

Exponential rate of convergence for some Markov operators

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The exponential rate of convergence for some Markov operators is established. The operators correspond to continuous iterated function systems which are a very useful tool in some cell cycle models.

I. INTRODUCTION

We are concerned with Markov operators corresponding to continuous iterated function systems. The main purpose of the paper is to prove spectral gap assuring exponential rate of convergence. The operators under consideration were used in Lasota & Mackey [9], where the authors studied some cell cycle model. See also Tyson & Hannsgen [16] or Murray & Hunt [11] to get more details on the subject. Lasota and Mackey proved only stability, while we managed to evaluate rate of convergence, bringing some information important from biological point of view. In our paper we base on coupling methods introduced in Hairer [4]. In the same spirit, exponential rate of convergence was proved in Ślęczka [15] for classical iterated function systems (see also Hairer & Mattingly [5] or Kapica & Ślęczka [7]). It is worth mentioning here that our result will allow us to show the Central Limit Theorem (CLT) and the Law of Iterated Logarithm (LIL). To do this, we will adapt general results recently proved in Bołt, Majewski & Szarek [2] or in Komorowski & Walczuk [8]. The proof of CLT and LIL will be provided in a future paper.

The organization of the paper goes as follows. Section 2 introduces basic notation and definitions that are needed throughout the paper. Most of them are adapted from Billingsley [1], Meyn & Tweedie [12], Lasota & Yorke [10] and Szarek [14]. Biological background is shortly presented in Section 3. Sections 4 and 5 provide the mathematical derivation of the model and the main theorem (Theorem 2), which establishes the exponential rate of convergence in the model. Sections 6-8 are devoted to the construction of coupling measure for iterated function systems. Thanks to the results presented in Section 9 we are finally able to present the proof of the main theorem in Section 10.

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II. NOTATION AND BASIC DEFINITIONS

Let (X, ϱ) be a Polish space. We denote by B_X the family of all Borel subsets of X . Let $C(X)$ be the space of all bounded and continuous functions $f : X \rightarrow \mathbb{R}$ with the supremum norm.

We denote by $M(X)$ the family of all Borel measures on X and by $M_{fin}(X)$ and $M_1(X)$ its subfamilies such that $\mu(X) < \infty$ and $\mu(X) = 1$, respectively. Elements of $M_{fin}(X)$ which satisfy $\mu(X) \leq 1$ are called sub-probability measures. To simplify notation, we write

$$\langle f, \mu \rangle = \int_X f(x) \mu(dx) \quad \text{for } f \in C(X), \mu \in M(X).$$

An operator $P : M_{fin}(X) \rightarrow M_{fin}(X)$ is called a Markov operator if

1. $P(\lambda_1 \mu_1 + \lambda_2 \mu_2) = \lambda_1 P\mu_1 + \lambda_2 P\mu_2$ for $\lambda_1, \lambda_2 \geq 0$, $\mu_1, \mu_2 \in M_{fin}(X)$;
2. $P\mu(X) = \mu(X)$ for $\mu \in M_{fin}(X)$.

If, additionally, there exists a linear operator $U : C(X) \rightarrow C(X)$ such that

$$\langle Uf, \mu \rangle = \langle f, P\mu \rangle \quad \text{for } f \in C(X), \mu \in M_{fin}(X),$$

an operator P is called a Feller operator. Every Markov operator P may be extended to the space of signed measures on X denoted by $M_{sig}(X) = \{\mu_1 - \mu_2 : \mu_1, \mu_2 \in M_{fin}(X)\}$. For $\mu \in M_{sig}(X)$ we denote by $\|\mu\|$ the total variation norm of μ , i.e.

$$\|\mu\| = \mu^+(X) + \mu^-(X),$$

where μ^+ and μ^- come from the Hahn-Jordan decomposition of μ (see Halmos [6]). For fixed $\bar{x} \in X$ we also consider the space $M_1^1(X)$ of all probability measures with the first moment finite, i.e. $M_1^1(X) = \{\mu \in M_1(X) : \int_X \varrho(x, \bar{x}) \mu(dx) < \infty\}$. The family is independent of the choice of $\bar{x} \in X$. We call $\mu_* \in M_{fin}(X)$ an invariant measure of P if $P\mu_* = \mu_*$. For $\mu \in M_{fin}(X)$ we define the support of μ by

$$\text{supp } \mu = \{x \in X : \mu(B(x, r)) > 0 \text{ for } r > 0\},$$

where $B(x, r)$ is the open ball in X with center at $x \in X$ and radius $r > 0$.

In $M_{sig}(X)$ we introduce the Fournet-Mourier norm

$$\|\mu\|_{\mathcal{L}} = \sup_{f \in \mathcal{L}} |\langle f, \mu \rangle|,$$

where

$$\mathcal{L} = \{f \in C(X) : |f(x) - f(y)| \leq \varrho(x, y), |f(x)| \leq 1 \text{ for } x, y \in X\}. \quad (1)$$

The space $M_1(X)$ with the metric $\|\mu_1 - \mu_2\|_{\mathcal{L}}$ is complete (see Fortet & Mourier [3] or Rachev [13]).

III. SHORTLY ABOUT THE MODEL OF CELL DIVISION CYCLE

Let $(\Omega, \mathcal{F}, \text{Prob})$ be a probability space. Suppose that each cell in a considered population consists of d different substances, whose masses are described by the vector $y(t) = (y^1(t), \dots, y^d(t))$, where $t \in [0, T]$ denotes an age of a cell. We assume that the evolution of the vector $y(t)$ is given by the formula $y(t) = \Pi(x, t)$, where $\Pi(x, 0) = x$. Here $\Pi : X \times [0, T) \rightarrow X$ is a given function. A simple example fulfilling these criteria is given by assuming that $y(t)$ satisfies a system of ordinary differential equations

$$\frac{dy}{dt} = g(t, y) \quad (2)$$

with the initial condition $y(0) = x$ and the solution of (2) is given by $y(t) = \Pi(x, t)$.

If x_n denotes the initial value $x = y(0)$ of substances in the n -th generation and t_n denotes the mitotic time in the n -th generation, the distribution is given by

$$\text{Prob}(t_n \in I | x_n = x) = \int_I p(x, s) ds \quad \text{for } I \in [0, T], n \in N. \quad (3)$$

The vector $y(t_n) = \Pi(x_n, t_n)$ with $y(0) = \Pi(x, 0) = x$ describes an amount of intercellular substance just before cell division in the n -th generation. We assume that each daughter cell contains exactly half of the components of its stem cell. Hence

$$x_{n+1} = \frac{1}{2} \Pi(x_n, t_n) \quad \text{for } n = 0, 1, 2, \dots \quad (4)$$

The behaviour of (3) and (4) may be also described by the sequence $(\mu_n)_{n \geq 1}$ of distributions

$$\mu_n(A) = \text{Prob}(x_n \in A) \quad \text{for } n = 0, 1, 2, \dots \text{ and } A \in B_X.$$

See Lasota & Mackey [9] for more details.

IV. ASSUMPTIONS

We assume that (X, ϱ) is a Polish space. Fix $T < \infty$. We consider a family $\{t_n : n = 0, 1, \dots\}$ of independent random variables taking values in $[0, T]$. The family is defined on the probability space $(\Omega, \mathcal{F}, \text{Prob})$. Note that $\text{Prob}(t_n < T | x_n = x) = 1$. Let $S : X \times [0, T) \rightarrow X$ be a continuous function and

$$x_{n+1} = S(x_n, t_n), \quad n = 0, 1, 2, \dots$$

We assume that $p : X \times [0, T) \rightarrow [0, \infty)$ is a lower semi-continuous, non-negative function such that, for every $x \in X$, $p(x, 0) = 0$ and $p(x, t) > 0$ for $t > 0$. In addition, p is normalized, i.e. $\int_0^T p(x, u) du = 1$ for $x \in X$. Let us further assume that for each $A \in B_X$

$$\text{Prob}(x_{n+1} \in A) := \mu_{n+1}(A), \quad \text{and} \quad P\mu_n = \mu_{n+1},$$

where

$$P\mu(A) = \int_X \left(\int_0^T 1_A(S(x, t)) p(x, t) dt \right) \mu(dx). \quad (5)$$

The following assumptions will be needed throughout the paper:

- (I) $\varrho(S(x, t), S(y, t)) \leq \lambda(t)\varrho(x, y)$ for $x, y \in X$, where $\lambda : [0, T) \rightarrow [0, \infty)$ is a Borel measurable function;
- (II) $a := \sup_{x \in X} \int_0^T \lambda(t) p(x, t) dt < 1$;
- (III) $\sup_{t \in [0, T)} \varrho(S(\bar{x}, t), \bar{x}) < \infty$ for some $\bar{x} \in X$;
- (IV) there exists σ such that $p : X \times [0, T) \rightarrow [\sigma, \infty)$ is a continuous function and $\bar{c} > 0$ such that $\int_0^T |p(x, t) - p(y, t)| dt \leq \bar{c}\varrho(x, y)$ for $x, y \in X$;
- (V) function p is bounded and we assume that $\delta = \inf\{p(x, t) : x \in X, t \in (0, T)\} > 0$, $M = \sup\{p(x, t) : x \in X, t \in (0, T)\}$.

V. MAIN THEOREM

Let P be the Markov operator in the cell division model defined above. Lasota and Mackey proved asymptotic stability of P , i.e. the existence of an invariant measure $\mu_* \in M_1(X)$ and weak convergence of $(P^n \mu)$ to μ_* for $\mu \in M_1(X)$. The theorem says.

Theorem 1. *Let $S : X \times [0, T] \rightarrow X$ and $p : X \times [0, T] \rightarrow [0, \infty)$ satisfy the following conditions*

1. $\varrho(S(x, t), S(y, t)) \leq \lambda_0(x, t)\varrho(x, y)$ for $x, y \in X$, $t \in [0, T]$ and λ_0 and S related to p by the conditions $\int_0^T \lambda_0(x, t) p(x, t) dt \leq r_0$ and $\int_0^T |S(0, t)| p(x, t) dt \leq r_1$ for $x \in X$;
2. $\int_0^T |p(x, t) - p(y, t)| dt \leq r_2 \varrho(x, y)$ for $x, y \in X$;
3. for every $x \in X$ there exists a minimal division time $\tau_x \in [0, T]$ such that $p(x, t) = 0$ for $0 \leq t \leq \tau_x$ and $p(x, t) > 0$ for $\tau_x < t \leq T$.

We assume moreover that $r_0 < 1$ and $r_1, r_2 < \infty$. Then, the system (3) and (4) is asymptotically stable.

Obviously, conditions (i) and (ii) of Theorem 1 are satisfied by assumptions (I)-(IV) of the model in consideration. Note that condition (iii) is also fulfilled with $\tau_x = 0$, as for every $x \in X$ we have $p(x, 0) = 0$ and $p(x, t) > 0$ for every $t > 0$ and $x \in X$. That is why we can assume the existence of an invariant measure in the model.

Our aim is to show that rate of convergence is exponential.

Theorem 2. *Let $\mu \in M_1^1$. Under assumptions (I)-(V) there exist $C = C(\mu) > 0$ and $q \in [0, 1)$ such that*

$$\|P^n \mu - \mu_*\|_{\mathcal{L}} \leq Cq^n \quad \text{for } n \in N.$$

VI. MEASURES ON THE PATHSPACE AND COUPLING

We consider a family of measures $\{Q_x : x \in X\}$ on X . We assume measurability of the mappings $x \mapsto Q_x(A)$ for each $A \in B_X$. Fix $n, m \in N$. Now, suppose that $\{Q_x : x \in X\}$ is a family of measures on X^n and $\{R_x : x \in X\}$ is a family of measures on X^m . We can define a family of measures $\{(RQ)_x : x \in X\}$ on $X^n \times X^m$

$$(RQ)_x(A \times B) = \int_A R_{z_n}(B) Q_x(dz), \quad (6)$$

where $z = (z_1, \dots, z_n)$ and $A \in B_{X^n}$, $B \in B_{X^m}$.

We consider a family of sub-probability measures $\{P_x : x \in X\}$ on X . We assume that the mapping $x \mapsto P_x(A)$ is measurable for each $A \in B_X$. Furthermore, if each P_x is a probability measure, $\{P_x : x \in X\}$ is a transition probability function. Thus $P_x(A)$ is the probability of transition from x to A . We want to define a family of measures on X^∞ . Set $x \in X$. One-dimensional distributions $\{P_x^n : n \in N\}$ are defined by induction on n

$$P_x^0(A) = \delta_x(A), \dots, P_x^{n+1}(A) = \int_X P_z(A) P_x^n(dz), \quad (7)$$

where $A \in B_X$. Following (6), we easily obtain two and higher-dimensional distributions. Finally, we get the family $\{P_x^\infty : x \in X\}$ of sub-probability measures on X^∞ . This construction was motivated by Hairer [4]. The existence of measures P_x^∞ is established by the Kolmogorov theorem. More precisely, there exists some probability space, on which we can define a stochastic process ξ with distribution ϕ_ξ such that

$$\phi_\xi(A) = \text{Prob}(\xi^{-1}(A)) := P_x^\infty(A) \quad \text{for } A \in B_{X^\infty}.$$

Therefore, P_x^∞ is the distribution of the Markov chain ξ on X^∞ with transition probability function $\{P_x : x \in X\}$ and $\phi_{\xi_0} = \delta_x$ for $x \in X$. If an initial distribution is given by any $\mu \in M_{fin}(X)$, not necessarily by δ_x , we define

$$P_\mu^\infty(A) = \int_X P_x^\infty(A) \mu(dx) \quad \text{for } A \in B_{X^\infty}.$$

Definition 3. Let a transition probability function $\{P_x : x \in X\}$ be given. A family of probability measures $\{C_{x,y} : x, y \in X\}$ on $X \times X$ such that

- $C_{x,y}(A \times X) = P_x(A)$ for $A \in B_X$,
- $C_{x,y}(X \times B) = P_y(B)$ for $B \in B_X$,

where $x, y \in X$, is called coupling.

VII. ITERATED FUNCTION SYSTEMS

We consider a continuous mapping $S : X \times [0, T) \rightarrow X$ and a lower semi-continuous, non-negative normalized function $p : X \times [0, T) \rightarrow [0, \infty)$. For each $A \in B_X$ we build a transition operator $P_x(A) = \Pi(x, A)$. Since $P\mu$ is given by (5) and $(P\mu)(A) = \int_X P_x(A) \mu(dx)$, we define P_x to be

$$P_x(A) = \int_0^T 1_A(S(x, t)) p(x, t) dt = \int_0^T \delta_{S(x, t)}(A) p(x, t) dt.$$

Once again, we apply (6) and (7) to construct measures on products. As previously, P_μ^∞ exists for $\mu \in M_{fin}(X)$. Obviously, $P^n \mu$ is the n -th marginal of P_μ^∞ .

Fix $\bar{x} \in X$. We define $V : X \rightarrow [0, \infty)$ to be

$$V(x) = \varrho(x, \bar{x}).$$

Let us evaluate an integral $\langle V, P\mu \rangle = \int_X \varrho(x, \bar{x}) P\mu(dx) = \int_X U\varrho(x, \bar{x}) \mu(dx)$, where U is a dual operator to P . Since P is a Feller operator given by (5), we can define $U : C(X) \rightarrow C(X)$ by

$$Uf(x) = \int_0^T f(S(x, t)) p(x, t) dt.$$

Hence, from initial assumptions (I) and (II), we obtain

$$\begin{aligned}
\langle V, P\mu \rangle &= \int_X \left(\int_0^T \varrho(S(x, t), \bar{x}) p(x, t) dt \right) \mu(dx) \\
&\leq \int_X \left(\int_0^T (\varrho(S(x, t), S(\bar{x}, t)) + \varrho(S(\bar{x}, t), \bar{x})) p(x, t) dt \right) \mu(dx) \\
&\leq \int_X \left(\int_0^T \lambda(t) \varrho(x, \bar{x}) p(x, t) dt + \int_0^T \varrho(S(\bar{x}, t), \bar{x}) p(x, t) dt \right) \mu(dx) \\
&\leq a \int_X \varrho(x, \bar{x}) \mu(dx) + \int_X \tilde{c} \mu(dx) \\
&= a \langle V, \mu \rangle + c,
\end{aligned}$$

where $c = \int_X \tilde{c} \mu(dx)$ and $\tilde{c} = \sup_{t \in [0, T]} \varrho(S(\bar{x}, t), \bar{x})$ exists from assumption (III). Fix probability measures $\mu, \nu \in M_1^1(X)$ and Borel sets $A, B \in B_X$. We consider $b \in M_1(X^2)$ such that

$$b(A \times X) = \mu(A), \quad b(X \times B) = \nu(B)$$

and $\bar{b} \in M_1(X^2)$ such that

$$\bar{b}(A \times X) = P\mu(A), \quad \bar{b}(X \times B) = P\nu(B).$$

Furthermore, we define $\bar{V} : X^2 \rightarrow [0, \infty)$

$$\bar{V}(x, y) = V(x) + V(y) \quad \text{for } x, y \in X.$$

Note that

$$\langle \bar{V}, \bar{b} \rangle \leq a \langle \bar{V}, b \rangle + 2c. \quad (8)$$

For measures $b \in M_{fin}^1(X^2)$ finite on X^2 and with the first moment finite we define the linear functional

$$\phi(b) = \int_{X^2} \varrho(x, y) b(dx, dy).$$

Following the above definitions, we easily obtain

$$\phi(b) \leq \langle \bar{V}, b \rangle. \quad (9)$$

VIII. COUPLING FOR ITERATED FUNCTION SYSTEMS

On X^∞ we define the transition sub-probability function

$$Q_{x,y}(A \times B) = \int_0^T \min\{p(x, t), p(y, t)\} \delta_{(S(x, t), S(y, t))}(A \times B) dt \quad \text{for } A, B \in B_X. \quad (10)$$

It is easy to check that

$$Q_{x,y}(A \times X) \leq \int_0^T p(x,t) \delta_{S(x,t)}(A) dt = \int_0^T 1_A(S(x,t)) p(x,t) dt = P_x(A)$$

and analogously

$$Q_{x,y}(X \times B) \leq P_y(B).$$

Let Q_b denote the measure

$$Q_b(A \times B) = \int_{X^2} Q_{x,y}(A \times B) b(dx, dy) \quad \text{for } A, B \in B_X. \quad (11)$$

Note that for every $A, B \in B_X$ we obtain

$$\begin{aligned} Q_b^{n+1}(A \times B) &= \int_{X^2} Q_{x,y}^{n+1}(A \times B) b(dx, dy) \\ &= \int_{X^2} \int_{X^2} Q_{z_1, z_2}(A \times B) Q_{x,y}^n(dz_1, dz_2) b(dx, dy) \\ &= \int_{X^2} Q_{z_1, z_2}(A \times B) \int_{X^2} Q_{x,y}^n(dz_1, dz_2) b(dx, dy) \\ &= \int_{X^2} Q_{z_1, z_2}(A \times B) Q_b^n(dx, dy) = Q_{Q_b^n}(A \times B). \end{aligned}$$

Again, we are able to construct measures on products, as well as we are able to construct Q_b^∞ on X^∞ . Now, we check that

$$\phi(Q_b) \leq a\phi(b). \quad (12)$$

Let us observe that

$$\begin{aligned} \phi(Q_b) &= \int_{X^2} \int_{X^2} \varrho(x, y) Q_{u,v}(dx, dy) b(du, dv) \\ &= \int_{X^2} \int_0^T \left(\int_{X^2} \varrho(x, y) \min\{p(u, t), p(v, t)\} \delta_{(S(u,t), S(v,t))}(dx, dy) \right) dt b(du, dv) \\ &\leq \int_{X^2} \int_0^T (\varrho(S(u, t), S(v, t)) p(u, t)) dt b(du, dv) \\ &\leq \int_{X^2} \int_0^T \lambda(t) \varrho(u, v) p(u, t) dt b(du, dv) \\ &\leq a \int_{X^2} \varrho(u, v) b(du, dv) \\ &= a\phi(b). \end{aligned}$$

We can find such a measure $R_{x,y}$ that the sum of $Q_{x,y}$ and $R_{x,y}$ gives a new coupling measure $C_{x,y}$, i.e. $C_{x,y}(A \times X) = P_x(A)$ and $C_{x,y}(X \times B) = P_y(B)$ for $A, B \in B_X$.

Lemma 4. *There exists the family $\{R_{x,y} : x, y \in X\}$ of measures on X^2 such that we can define*

$$C_{x,y} = Q_{x,y} + R_{x,y} \quad \text{for } x, y \in X$$

and, moreover,

- (i) *the mapping $(x, y) \mapsto R_{x,y}(A \times B)$ is measurable for every $A, B \in B_X$;*
- (ii) *measures $R_{x,y}$ are non-negative for $x, y \in X$;*
- (iii) *measures $C_{x,y}$ are probabilistic for every $x, y \in X$ and so $\{C_{x,y} : x, y \in X\}$ is the transition probability function on X^2 ;*
- (iv) *for every $A, B \in B_X$ and $x, y \in X$ we get $C_{x,y}(A \times X) = P_x(A)$ and $C_{x,y}(X \times B) = P_y(B)$.*

Proof. Fix $A, B \in B_X$. Let

$$R_{x,y}(A \times B) = \begin{cases} (1 - Q_{x,y}(X^2))^{-1}(P_x(A) - Q_{x,y}(A \times X))(P_y(B) - Q_{x,y}(X \times B)), & Q_{x,y}(X^2) < 1 \\ 0, & Q_{x,y}(X^2) = 1. \end{cases}$$

Obviously, the formula may be extended to the measure. The mapping has all desirable properties (i)-(iv). \square

Lemma 4 shows that we can construct the coupling $\{C_{x,y} : x, y \in X\}$ for $\{P_x : x \in X\}$ such that $Q_{x,y} \leq C_{x,y}$, whereas measures $R_{x,y}$ are non-negative. By (6) and (7) we obtain the family of probability measures $\{C_{x,y}^\infty : x, y \in X\}$ on $(X^2)^\infty$ with marginals P_x^∞ and P_y^∞ . This construction appears in Hairer [4].

Fix $(x_0, y_0) \in X^2$. The transition probability function $\{C_{x,y} : x, y \in X\}$ defines the Markov chain Φ on X^2 with starting point (x_0, y_0) , while the transition probability function $\{\hat{C}_{x,y,\theta} : x, y \in X, \theta \in \{0, 1\}\}$ defines the Markov chain $\hat{\Phi}$ on the augmented space $X^2 \times \{0, 1\}$ with initial distribution $\hat{C}_{x_0, y_0}^0 = \delta_{(x_0, y_0, 1)}$. If $\hat{\Phi}_n = (x, y, i)$, where $x, y \in X, i \in \{0, 1\}$, then

$$\text{Prob}(\hat{\Phi}_{n+1} \in A \times B \times \{1\} \mid \hat{\Phi}_n = (x, y, i), i \in \{0, 1\}) = Q_{x,y}(A \times B),$$

$$\text{Prob}(\hat{\Phi}_{n+1} \in A \times B \times \{0\} \mid \hat{\Phi}_n = (x, y, i), i \in \{0, 1\}) = R_{x,y}(A \times B),$$

where $A, B \in B_X$. Once again, we refer to (6) and (7) to obtain the measure $\hat{C}_{x_0, y_0}^\infty$ on $(X^2 \times \{0, 1\})^\infty$ which is associated with the Markov chain $\hat{\Phi}$.

From now on, we assume that processes Φ and $\hat{\Phi}$ taking values in X^2 and $X^2 \times \{0, 1\}$, respectively, are defined on $(\Omega, \mathcal{F}, \mathbf{P})$. The expected value of the measures C_{x_0, y_0}^∞ or $\hat{C}_{x_0, y_0}^\infty$ is denoted by E_{x_0, y_0} .

IX. AUXILIARY THEOREMS

Fix $\varepsilon \in (0, 1 - a)$. Set

$$K_\varepsilon = \{(x, y) \in X^2 : \bar{V}(x, y) < \varepsilon^{-1}2c\},$$

where c is defined in Section VII. Let $d : (X^2)^\infty \rightarrow N$ denote the time of the first visit in K_ε , i.e.

$$d((x_n, y_n)_{n \in N}) = \inf\{n \geq 1 : (x_n, y_n) \in K_\varepsilon\}.$$

Theorem 5. *For every $\gamma \in (0, 1)$ there exist positive constants C_1, C_2 such that*

$$E_{x_0, y_0} \left((a + \varepsilon)^{-\gamma d} \right) \leq C_1 \bar{V}(x_0, y_0) + C_2.$$

Proof. Fix $(x_0, y_0) \in X^2$. Let $\Phi = (X_n, Y_n)_{n \in N}$ be the Markov chain with starting point (x_0, y_0) and transition probability function $\{C_{x, y} : x, y \in X\}$. Let $F_n \subset F$, $n \in N$ be the natural filtration in Ω associated with Φ . We define

$$A_n = \{\omega \in \Omega : \Phi_i = (X_i(\omega), Y_i(\omega)) \notin K_\varepsilon \text{ for } i = 1, \dots, n\}, \quad n \in N.$$

Obviously $A_{n+1} \subset A_n$ and $A_n \in F_n$ for $n \in N$. The following inequalities are \mathbf{P} -a.s. satisfied in Ω

$$1_{A_n} E_{x_0, y_0} (\bar{V}(X_{n+1}, Y_{n+1}) | F_n) \leq 1_{A_n} (a \bar{V}(X_n, Y_n) + 2c) \leq 1_{A_n} (a + \varepsilon) \bar{V}(X_n, Y_n).$$

The first inequality is a consequence of (8), the second follows directly from the definitions of A_n and K_ε . Accordingly, we obtain

$$\begin{aligned} \int_{A_n} \bar{V}(X_n, Y_n) d\mathbf{P} &\leq \int_{A_{n-1}} \bar{V}(X_n, Y_n) d\mathbf{P} = \int_{A_{n-1}} E(\bar{V}(X_n, Y_n) | F_{n-1}) d\mathbf{P} \\ &\leq \int_{A_{n-1}} (a \bar{V}(X_{n-1}, Y_{n-1}) + 2c) d\mathbf{P} \leq (a + \varepsilon) \int_{A_{n-1}} \bar{V}(X_{n-1}, Y_{n-1}) d\mathbf{P}. \end{aligned}$$

On applying this estimates finitely many times, we obtain

$$\int_{A_n} \bar{V}(X_n, Y_n) d\mathbf{P} \leq (a + \varepsilon)^{n-1} \int_{A_1} \bar{V}(X_1, Y_1) d\mathbf{P} \leq (a + \varepsilon)^{n-1} (a \bar{V}(X_0, Y_0) + 2c).$$

Note that

$$\mathbf{P}(A_n) \leq \int_{A_n} \varepsilon (2c)^{-1} \bar{V}(X_n, Y_n) d\mathbf{P} \leq \varepsilon (2c(a + \varepsilon))^{-1} (a + \varepsilon)^n (a \bar{V}(X_0, Y_0) + 2c).$$

Set $\hat{c} := \varepsilon (2c(a + \varepsilon))^{-1} (a \bar{V}(X_0, Y_0) + 2c)$. Then, $\mathbf{P}(A_n) \leq (a + \varepsilon)^n \hat{c}$. Fix $\gamma \in (0, 1)$. Since d takes natural values $n \in N$, we obtain

$$\sum_{n=1}^{\infty} (a + \varepsilon)^{-\gamma n} \mathbf{P}(A_n) \leq \sum_{n=1}^{\infty} (a + \varepsilon)^{-\gamma n} (a + \varepsilon)^n \hat{c} = \sum_{n=1}^{\infty} (a + \varepsilon)^{(1-\gamma)n} \hat{c},$$

which implies convergence of the series. The proof is complete by the definition of \hat{c} and with properly choosen C_1, C_2 . \square

For every positive $r > 0$ we determine the set

$$C_r = \{(x, y) \in X^2 : \varrho(x, y) < r\}.$$

Lemma 6. Fix $\tilde{a} \in (a, 1)$. Let C_r be the set defined above and suppose that $\text{supp } b \subset C_r$. There exists $\bar{\gamma} > 0$ such that

$$Q_b(C_{\tilde{a}r}) \geq \bar{\gamma} \|b\|$$

for a , δ and M defined in Section IV.

Proof. Directly from (11) and (10) we obtain

$$\begin{aligned} Q_b(C_{\tilde{a}r}) &= \int_{X^2} \int_0^T \min\{p(x, t), p(y, t)\} \delta_{(S(x, t), S(y, t))}(C_{\tilde{a}r}) dt b(dx, dy) \\ &= \int_{X^2} \left(\int_0^T \min\{p(x, t), p(y, t)\} 1_{C_{\tilde{a}r}}(S(x, t), S(y, t)) dt \right) b(dx, dy). \end{aligned}$$

Note that $1_{C_{\tilde{a}r}}(S(x, t), S(y, t)) = 1$ if and only if $t \in \mathcal{T}$, where

$$\mathcal{T} := \{t \in (0, T) : \varrho(S(x, t), S(y, t)) < \tilde{a}r\}.$$

Set $\mathcal{T}' := (0, T) \setminus \mathcal{T}$. Hence

$$\begin{aligned} Q_b(C_{\tilde{a}r}) &= \int_{X^2} \left(\int_{\mathcal{T}} \min\{p(x, t), p(y, t)\} 1_{C_{\tilde{a}r}}(S(x, t), S(y, t)) dt \right. \\ &\quad \left. + \int_{\mathcal{T}'} \min\{p(x, t), p(y, t)\} 1_{C_{\tilde{a}r}}(S(x, t), S(y, t)) dt \right) b(dx, dy). \end{aligned}$$

Note that

$$\int_{\mathcal{T}'} \min\{p(x, t), p(y, t)\} \varrho(S(x, t), S(y, t)) dt \leq \int_{\mathcal{T}'} p(x, t) \lambda(t) \varrho(x, y) dt \leq a \varrho(x, y),$$

so for $(x, y) \in C_r$

$$\int_{\mathcal{T}'} \min\{p(x, t), p(y, t)\} \varrho(S(x, t), S(y, t)) dt \leq ar.$$

However,

$$\tilde{a}r \int_{\mathcal{T}'} p(x, t) dt < \int_{\mathcal{T}'} p(x, t) \varrho(S(x, t), S(y, t)) dt.$$

Therefore

$$\int_{\mathcal{T}'} p(x, t) dt < \frac{a}{\tilde{a}} < 1,$$

which implies that the first integral is non-zero. Furthermore, the length of \mathcal{T}' satisfies $|\mathcal{T}'| < a(\tilde{a}\delta)^{-1}$. We obtain

$$\int_{\mathcal{T}} p(x, t) dt \geq 1 - \frac{a}{\tilde{a}} = \gamma,$$

which means that $|\mathcal{T}| \geq M^{-1}\gamma$. Finally,

$$\begin{aligned} Q_b(C_{\tilde{a}r}) &\geq \int_{X^2} \int_{\mathcal{T}} \min\{p(x, t), p(y, t)\} 1_{C_{\tilde{a}r}}(S(x, t), S(y, t)) dt b(dx, dy) \\ &\geq \int_{X^2} \delta |\mathcal{T}| b(dx, dy) \geq \delta \frac{\gamma}{M} \|b\|. \end{aligned}$$

If we set $\bar{\gamma} := \delta M^{-1}\gamma$, the proof is complete. \square

Theorem 7. *For every $\varepsilon \in (0, 1 - a)$ there exists $n_0 \in \mathbb{N}$ such that*

$$\|Q_{x,y}^\infty\| \geq \frac{1}{2} \bar{\gamma}^{n_0} \quad \text{for } (x, y) \in K_\varepsilon,$$

where $\bar{\gamma} > 0$ is given in Lemma 6.

Proof. Note that for every $(x, y) \in X^2$

$$\int_0^T (\min\{p(x, t), p(y, t)\} + |p(x, t) - p(y, t)| - p(x, t)) dt \geq 0,$$

and therefore

$$\|Q_{x,y}\| + \int_0^T |p(x, t) - p(y, t)| dt \geq 1.$$

From assumption (IV) there is $\bar{c} > 0$ such that

$$\|Q_{x,y}\| \geq 1 - \int_0^T |p(x, t) - p(y, t)| dt \geq 1 - \bar{c} \varrho(x, y).$$

For every $b \in M_{fin}(X^2)$ we get

$$\|Q_b\| = \int_{X^2} \|Q_{x,y}\| b(dx, dy) \geq \int_{X^2} b(dx, dy) - \bar{c} \int_{X^2} \varrho(x, y) b(dx, dy) = \|b\| - \bar{c} \phi(b).$$

Property (12) implies that

$$\|Q_b^{n+1}\| \geq \|b\| - \bar{c} \left(\sum_{k=0}^n a^k \right) \phi(b) \geq \|b\| - (1 - a)^{-1} \bar{c} \phi(b), \quad n \in \mathbb{N}.$$

If $\text{supp } b \subset C_r$, then

$$\phi(b) \leq \int_{C_r} \varrho(x, y) b(dx, dy) \leq r \|b\|.$$

Let $r = (2\bar{c})^{-1}(1 - a)$. We obtain

$$\|Q_b^\infty\| \geq \frac{\|b\|}{2}.$$

Fix $\varepsilon \in (0, 1 - a)$. It is clear that $K_\varepsilon \subset C_{\varepsilon^{-1}2c}$. If we define $n_0 := \min\{n \geq 1 : a^n(\varepsilon)^{-1}2c < r\}$, then $C_{a^{n_0}\varepsilon^{-1}2c} \subset C_r$. Remembering that $Q_{x,y}^{n+m} = Q_{Q_{x,y}^n}^m$ and using the Markov property, we obtain

$$Q_{x,y}^\infty(X^2) \geq Q_{Q_{x,y}^{n_0}}^\infty(X^2).$$

According to Lemma 6, we obtain

$$\|Q_{x,y}^\infty\| \geq \|Q_{Q_{x,y}^{n_0}}^\infty\| \geq \frac{\|Q_{x,y}^{n_0}\|}{2} = \frac{Q_{x,y}^{n_0}(C_r)}{2} \geq \frac{Q_{x,y}^{n_0}(C_{a^{n_0}\varepsilon^{-1}2c})}{2} \geq \frac{\bar{\gamma}^{n_0}}{2}$$

for $(x, y) \in K_\varepsilon$. This finishes the proof. \square

Definition 8. *Coupling time $\tau : (X^2 \times \{0, 1\})^\infty \rightarrow N$ is defined as follows*

$$\tau((x_n, y_n, \theta_n)_{n \in N}) = \inf\{n \geq 1 : \theta_k = 1 \text{ for } k \geq n\}.$$

Theorem 9. *There exist $\tilde{q} \in (0, 1)$ and $C_3 > 0$ such that*

$$E_{x,y}(\tilde{q}^{-\tau}) \leq C_3(1 + \bar{V}(x, y)) \text{ for } (x, y) \in X^2.$$

Proof. Fix $\varepsilon \in (0, 1 - a)$ and $(x, y) \in X$. To simplify notation, we write $\beta = (a + \varepsilon)^{\frac{1}{2}}$. Let d be the random moment of the first visit in K_ε . Suppose that

$$d_1 = d, \quad d_{n+1} = d_n + d \circ T_{d_n},$$

where $n > 1$ and T_n are shift operators on $(X^2 \times \{0, 1\})^\infty$, i.e. $T_n((x_k, y_k, \theta_k)_{k \in N}) = (x_{k+n}, y_{k+n}, \theta_{k+n})_{k \in N}$. Theorem 5 implies that every d_n is $C_{x,y}^\infty$ -a.s. finished. The strong Markov property shows that

$$E_{x,y}(\beta^d \circ T_{d_n} | F_{d_n}) = E_{(X_{d_n}, Y_{d_n})}(\beta^d) \text{ for } n \in N,$$

where F_{d_n} denotes the σ -algebra on $(X^2 \times \{0, 1\})$ generated by d_n and $\Phi = (X_n, Y_n)_{n \in N}$ is the Markov chain with transition probability function $\{C_{x,y}^\infty : x, y \in X\}$. By Theorem 5 and the definition of K_ε we obtain

$$E_{x,y}(\beta^{d_{n+1}}) = E_{x,y}(\beta^{d_n} E_{(X_{d_n}, Y_{d_n})}(\beta^d)) \leq E_{x,y}(\beta^{d_n}) (C_1 \varepsilon^{-1} 2c + C_2).$$

Fix $\eta = C_1 \varepsilon^{-1} 2c + C_2$. Consequently,

$$E_{x,y}(\beta^{d_{n+1}}) \leq \eta^n E_{x,y}(\beta^d) \leq \eta^n (C_1 \bar{V}(x, y) + C_2). \quad (13)$$

We define $\hat{\tau}((x_n, y_n, \theta_n)_{n \in N}) = \inf\{n \geq 1 : (x_n, y_n) \in K_\varepsilon, \theta_k = 1 \text{ for } k \geq n\}$ and $\sigma = \inf\{n \geq 1 : \hat{\tau} = d_n\}$. By Theorem 7 there is $n_0 \in N$ such that

$$\hat{C}_{x,y}^\infty(\sigma > n) \leq (1 - \frac{\bar{\gamma}^{n_0}}{2})^n \quad \text{for } n \in N. \quad (14)$$

Let $d > 1$. By the Hölder inequality, (13) and (14) we obtain

$$\begin{aligned} E_{x,y} \left(\beta^{\frac{\hat{\tau}}{p}} \right) &\leq \sum_{k=1}^{\infty} E_{x,y} \left(\beta^{\frac{d_k}{p}} 1_{\sigma=k} \right) \leq \sum_{k=1}^{\infty} \left(E_{x,y} \left(\beta^{d_k} \right) \right)^{\frac{1}{p}} \left(\hat{C}_{x,y}^\infty(\sigma = k) \right)^{(1-\frac{1}{p})} \\ &\leq (C_1 \bar{V}(x, y) + C_2)^{\frac{1}{p}} \eta^{-\frac{1}{p}} \sum_{k=1}^{\infty} \eta^{\frac{k}{p}} (1 - \frac{1}{2} \bar{\gamma}^{n_0})^{(k-1)(1-\frac{1}{p})} \\ &= (C_1 \bar{V}(x, y) + C_2)^{\frac{1}{p}} \eta^{-\frac{1}{p}} (1 - \frac{1}{2} \bar{\gamma}^{n_0})^{-(1-\frac{1}{p})} \sum_{k=1}^{\infty} \left(\left(\frac{\eta}{1 - \frac{1}{2} \bar{\gamma}^{n_0}} \right)^{\frac{1}{p}} (1 - \frac{1}{2} \bar{\gamma}^{n_0}) \right)^k. \end{aligned}$$

For p sufficiently large and $\tilde{q} = \beta^{-\frac{1}{p}}$, we get

$$E_{x,y} \left(\tilde{q}^{-\hat{\tau}} \right) = E_{x,y} \left(\beta^{\frac{\hat{\tau}}{p}} \right) \leq (1 + \bar{V}(x, y)) C_3$$

for some C_3 . Since $\tau \leq \hat{\tau}$, we finish the proof. \square

Theorem 10. *There exist $q \in (0, 1)$ and $C_6 > 0$ such that*

$$\|P_x^n - P_y^n\|_{\mathcal{L}} \leq q^n C_6 (1 + \bar{V}(x, y)) \quad \text{for } x, y \in X \text{ and } n \in N.$$

Proof. For $n \in N$ we define sets

$$A_{\frac{n}{2}} = \{t \in (X^2 \times \{0, 1\})^\infty : \tau(t) \leq \frac{n}{2}\},$$

$$B_{\frac{n}{2}} = \{t \in (X^2 \times \{0, 1\})^\infty : \tau(t) > \frac{n}{2}\}.$$

Note that $A_{\frac{n}{2}} \cap B_{\frac{n}{2}} = \emptyset$ and $A_{\frac{n}{2}} \cup B_{\frac{n}{2}} = (X^2 \times \{0, 1\})^\infty$, so for $n \in N$ we have

$$\hat{C}_{x,y}^\infty = \hat{C}_{x,y}^\infty|_{A_{\frac{n}{2}}} + \hat{C}_{x,y}^\infty|_{B_{\frac{n}{2}}}.$$

Hence,

$$\|P_x^n - P_y^n\|_{\mathcal{L}} = \sup_{f \in \mathcal{L}} \left| \int_{X^2} f(z) (P_x^n - P_y^n)(dz) \right| = \sup_{f \in \mathcal{L}} \left| \int_{X^2} (f(z_1) - f(z_2)) (\Pi_{X^2}^* \Pi_n^* \hat{C}_{x,y}^\infty)(dz_1, dz_2) \right|,$$

where $\Pi_n : (X^2 \times \{0, 1\})^\infty \rightarrow X^2 \times \{0, 1\}$ are the projections on the n -th component and $\Pi_{X^2} : X^2 \times \{0, 1\} \rightarrow X^2$ is the projection on X^2 . Now, recalling the definition of the set \mathcal{L} (see (1)), we

obtain

$$\begin{aligned}
\|P_x^n - P_y^n\|_{\mathcal{L}} &= \sup_{f \in \mathcal{L}} \left| \int_{X^2} (f(z_1) - f(z_2)) (\Pi_{X^2}^* \Pi_n^* \hat{C}_{x,y}^\infty|_{A_{\frac{n}{2}}}) (dz_1, dz_2) \right. \\
&\quad \left. + \int_{X^2} (f(z_1) - f(z_2)) (\Pi_{X^2}^* \Pi_n^* \hat{C}_{x,y}^\infty|_{B_{\frac{n}{2}}}) (dz_1, dz_2) \right| \\
&\leq \sup_{f \in \mathcal{L}} \left| \int_{X^2} (f(z_1) - f(z_2)) (\Pi_{X^2}^* \Pi_n^* \hat{C}_{x,y}^\infty|_{A_{\frac{n}{2}}}) (dz_1, dz_2) \right| + 2\hat{C}_{x,y}^\infty(B_{\frac{n}{2}}) \\
&\leq \sup_{f \in \mathcal{L}} \left| \int_{X^2} \varrho(z_1, z_2) (\Pi_{X^2}^* \Pi_n^* \hat{C}_{x,y}^\infty|_{A_{\frac{n}{2}}}) (dz_1, dz_2) \right| + 2\hat{C}_{x,y}^\infty(B_{\frac{n}{2}}).
\end{aligned}$$

Note that by iterative application of (12) we obtain

$$\int_{X^2} \varrho(z_1, z_2) (\Pi_{X^2}^* \Pi_n^* \hat{C}_{x,y}^\infty|_{A_{\frac{n}{2}}}) (dz_1, dz_2) = \phi(\Pi_{X^2}^* \Pi_n^* (\hat{C}_{x,y}^\infty|_{A_{\frac{n}{2}}})) \leq a^{\frac{n}{2}} \phi(\Pi_{X^2}^* \Pi_{\frac{n}{2}}^* (\hat{C}_{x,y}^\infty|_{A_{\frac{n}{2}}})).$$

Then it follows from (8) and (9) that

$$\phi(\Pi_{X^2}^* \Pi_{\frac{n}{2}}^* (\hat{C}_{x,y}^\infty|_{A_{\frac{n}{2}}})) \leq a^{\frac{n}{2}} \bar{V}(x, y) + \frac{2c}{1-a}$$

We obtain coupling inequality

$$\|P_x^n - P_y^n\|_{\mathcal{L}} \leq a^{\frac{n}{2}} \left(a^{\frac{n}{2}} \bar{V}(x, y) + \frac{2c}{1-a} \right) + 2\hat{C}_{x,y}^\infty(B_{\frac{n}{2}}).$$

It follows from Theorem 10 and the Chebyshev inequality that

$$\hat{C}_{x,y}^\infty(B_{\frac{n}{2}}) = \hat{C}_{x,y}^\infty(\{\tau > \frac{n}{2}\}) = \hat{C}_{x,y}^\infty(\{\tilde{q}^{-\tau} \leq \tilde{q}^{-\frac{n}{2}}\}) \leq \frac{E_{x,y}(\tilde{q}^{-\tau})}{\tilde{q}^{-\frac{n}{2}}} \leq \tilde{q}^{\frac{n}{2}} C_4 (1 + \bar{V}(x, y))$$

for some $\tilde{q} \in (0, 1)$ and $C_4 > 0$. Finally,

$$\|P_x^n - P_y^n\|_{\mathcal{L}} \leq a^{\frac{n}{2}} C_5 (1 + \bar{V}(x, y)) + 2\tilde{q}^{\frac{n}{2}} C_4 (1 + \bar{V}(x, y)),$$

where $C_5 = \max\{a^{\frac{n}{2}}, (1-a)^{-1}2c\}$. Setting $q := \max\{a^{\frac{1}{2}}, \tilde{q}^{\frac{1}{2}}\}$ and $C_6 := C_5 + 2C_4$, gives our claim. \square

X. PROOF OF THE MAIN THEOREM

Theorem 10 is essential to the following proof.

Proof. Theorem 10 implies that

$$\|P_x^n - P_y^n\|_{\mathcal{L}} \leq q^n C_6 (1 + \bar{V}(x, y)) \quad \text{for } x, y \in X \text{ and } n \in N,$$

where q and C_6 are the appropriate constants. Obviously,

$$\|P^n \mu - \mu_*\|_{\mathcal{L}} = \|P^n \mu - P^n \mu_*\|_{\mathcal{L}} = \sup_{f \in \mathcal{L}} \left| \int_X f(z) P^n \mu(dz) - \int_X f(z) P^n \mu_*(dz) \right|.$$

Note that

$$\begin{aligned}
\int_X f(z)P^n\mu(dz) - \int_X f(z)P^n\mu_*(dz) &= \int_X \int_X f(z)P_x^n(dz)\mu(dx) - \int_X \int_X f(z)P_y^n(dz)\mu_*(dy) \\
&= \int_X \int_X \left(\int_X f(z)P_x^n(dz) - \int_X f(z)P_y^n(dz) \right) \mu_*(dy)\mu(dx) \\
&\leq \int_X \int_X \|P_x^n - P_y^n\|_{\mathcal{L}} \mu_*(dy)\mu(dx) \\
&\leq q^n C,
\end{aligned}$$

where $C := \int_X \int_X C_6(1 + \bar{V}(x, y))\mu_*(dy)\mu(dx)$. Since C is dependant only on μ , the proof is complete. \square

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